# On the necessity of negative coefficients for operator splitting schemes of order higher than two 

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#### Abstract

In this paper we analyse numerical integration methods applied to differential equations which are separable in solvable parts. These methods are compositions of flows associated with each part of the system. We propose an elementary proof of the necessary existence of negative coefficients if the schemes are of order, or effective order, $p \geqslant 3$ and provide additional information about the distribution of these negative coefficients. It is shown that if the methods involve flows associated with more general terms this result does not necessarily apply and in some cases it is possible to build higher order schemes with positive coefficients. © 2004 IMACS. Published by Elsevier B.V. All rights reserved.


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## 1. Introduction

Operator splitting schemes are numerical methods which are particularly useful to approximate the evolution of differential equations when they are separable in solvable parts [17]. To be more specific, let us consider the ODE

$$
\begin{equation*}
x^{\prime}=f(x), \quad x_{0}=x(0) \in \mathbb{R}^{D} \tag{1}
\end{equation*}
$$

[^0]with $f: \mathbb{R}^{D} \rightarrow \mathbb{R}^{D}$ and associated vector field (or Lie operator associated with $f$ )
\[

$$
\begin{equation*}
F=\sum_{i=1}^{D} f_{i}(x) \frac{\partial}{\partial x_{i}} \tag{2}
\end{equation*}
$$

\]

Let us denote by $\varphi_{h}$ the $h$-flow of the ODE (1) for a given time step $h$. In other words, the exact solution is given by $x(h)=\varphi_{h}\left(x_{0}\right)$.

Now let us assume that $f(x)$ can be written as $f(x)=f_{A}(x)+f_{B}(x)$, the vector field $F$ is split accordingly as $F=F_{A}+F_{B}$ and the $h$-flows $\varphi_{h}^{[A]}$ and $\varphi_{h}^{[B]}$ corresponding to $F_{A}$ and $F_{B}$, respectively, can be exactly computed or, equivalently, the equations $x^{\prime}=f_{A}(x)$ and $x^{\prime}=f_{B}(x)$ are solvable. Then the composition (sometimes called Lie-Trotter splitting)

$$
\begin{equation*}
\psi_{h, 1} \equiv \varphi_{h}^{[B]} \circ \varphi_{h}^{[A]} \tag{3}
\end{equation*}
$$

approximates $\varphi_{h}$ with error of order $h^{2}$, i.e., $\psi_{h, 1}\left(x_{0}\right)=\varphi_{h}\left(x_{0}\right)+\mathcal{O}\left(h^{2}\right)$, whereas the so-called Strang splitting or Störmer/Verlet/leapfrog scheme

$$
\begin{equation*}
\psi_{h, 2} \equiv \varphi_{h / 2}^{[A]} \circ \varphi_{h}^{[B]} \circ \varphi_{h / 2}^{[A]} \tag{4}
\end{equation*}
$$

is such that $\psi_{h, 2}\left(x_{0}\right)=\varphi_{h}\left(x_{0}\right)+\mathcal{O}\left(h^{3}\right)$. The order of approximation to the exact solution can be increased by including more maps with fractional time steps in the composition. In general, the scheme

$$
\begin{equation*}
\psi_{h} \equiv \varphi_{b_{m} h}^{[B]} \circ \varphi_{a_{m} h}^{[A]} \circ \cdots \circ \varphi_{b_{1} h}^{[B]} \circ \varphi_{a_{1} h}^{[A]} \tag{5}
\end{equation*}
$$

$\left(a_{1}, b_{1}, \ldots, a_{m}, b_{m}\right) \in \mathbb{R}^{2 m}$, is of order $p$ if $\psi_{h}=\varphi_{h}+\mathcal{O}\left(h^{p+1}\right)$ for a proper choice of $m$ and coefficients $a_{i}, b_{i}$. It can be assumed without loss of generality that in (5) none of the coefficients $b_{1}, b_{2}, \ldots, b_{m-1}$ as well as none of $a_{2}, a_{3}, \ldots, a_{m}$ are vanishing, i.e., only $a_{1}$ and/or $b_{m}$ can be zero, since otherwise the corresponding flows could be removed $\left(\varphi_{0}^{[A]}=\varphi_{0}^{[B]}=\mathrm{id}\right.$, the identity map) and the rest of the maps would be concatenated (due to the group property of the flows).

Numerical schemes of order $p \geqslant 3$ based on the composition (5) have been successfully applied for solving a large number of problems [10,17], including also certain partial differential equations. In fact, splitting methods are frequently used in celestial mechanics [21], quantum mechanics [6], molecular dynamics [12], accelerator physics [7] and, in general, for numerically solving Hamiltonian dynamical systems [8,16], Poisson systems [14] and reversible differential equations [15]. It has been noticed, however, that some of the coefficients in (5) are negative for $p \geqslant 3$ when one considers arbitrary vector fields $F_{A}$ and $F_{B}$. In other words, the methods always involve stepping backwards in time. This constitutes a problem when the differential equation is defined in a semigroup, as arises sometimes in applications, since then the method can only be conditionally stable [17]. Also schemes with negative coefficients may not be well-posed when applied to PDEs involving unbounded operators.

The existence of backward fractional time steps in the composition method (5) is in fact unavoidable, and can be established as the following two theorems:

Theorem 1 [20,22]. If $p$ is a positive integer such that $p \geqslant 3$, then there are no composition methods of the form (5) and finite $m$ with all the coefficients $a_{i}, b_{i}$ being positive.

Theorem 2 [9]. If $p$ is a positive integer such that $p \geqslant 3$, then, for every pth-order method (5) with $m$ any finite positive integer, one has

$$
\min _{1 \leqslant i \leqslant m} a_{i}<0 \quad \text { and } \quad \min _{1 \leqslant j \leqslant m} b_{j}<0 .
$$

Theorem 2 is stronger than Theorem 1, in the sense that it establishes that at least one of the $a_{i}$ and also one of the $b_{i}$ coefficients have to be strictly negative, although a similar (and certainly non-trivial) proof strategy was used. One of the goals of this paper is to provide an alternative, elementary proof of Theorem 2, giving in addition a more detailed analysis of how negative coefficients are distributed in the composition.

During the last few years the processing technique has been used to find composition methods requiring less evaluations than conventional schemes of order $p$ [13]. The idea consists in enhancing an integrator $\psi_{h}$ (the kernel) with a parametric map $\pi_{h}: \mathbb{R}^{D} \rightarrow \mathbb{R}^{D}$ (the post-processor) as

$$
\begin{equation*}
\hat{\psi}_{h}=\pi_{h} \circ \psi_{h} \circ \pi_{h}^{-1} \tag{6}
\end{equation*}
$$

Application of $n$ steps of the new (and hopefully better) integrator $\hat{\psi}_{h}$ leads to

$$
\hat{\psi}_{h}^{n}=\pi_{h} \circ \psi_{h}^{n} \circ \pi_{h}^{-1},
$$

which can be considered as a change of coordinates in phase space. Observe that processing is advantageous if $\hat{\psi}_{h}$ is a more accurate method than $\psi_{h}$ and the cost of $\pi_{h}$ is negligible, since it provides the accuracy of $\hat{\psi}_{h}$ at essentially the cost of the least accurate method $\psi_{h}$.

The simplest example of a processed integrator is provided in fact by the Strang splitting (4). As a consequence of the group property of the exact flow, we have

$$
\begin{align*}
\psi_{h, 2} & =\varphi_{h / 2}^{[A]} \circ \varphi_{h}^{[B]} \circ \varphi_{h / 2}^{[A]}=\varphi_{h / 2}^{[A]} \circ \varphi_{h}^{[B]} \circ \varphi_{h}^{[A]} \circ \varphi_{-h}^{[A]} \circ \varphi_{h / 2}^{[A]} \\
& =\varphi_{h / 2}^{[A]} \circ \psi_{h, 1} \circ \varphi_{-h / 2}^{[A]}=\pi_{h} \circ \psi_{h, 1} \circ \pi_{h}^{-1} \tag{7}
\end{align*}
$$

with $\pi_{h}=\varphi_{h / 2}^{[A]}$. Hence, applying the first order method (3) with processing yields a second order of approximation.

Although initially intended for Runge-Kutta methods [4], the processing technique has proved its usefulness mainly in the context of geometric numerical integration [10], where constant step-sizes are widely employed.

We say that the method $\psi_{h}$ is of effective order $p$ if a post-processor $\pi_{h}$ exists for which $\hat{\psi}_{h}$ is of (conventional) order $p$ [4], that is,

$$
\pi_{h} \circ \psi_{h} \circ \pi_{h}^{-1}=\varphi_{h}+\mathcal{O}\left(h^{p+1}\right)
$$

Hence, as the previous example shows, the Lie-Trotter splitting (3) is of effective order 2. Obviously, a method of order $p$ is also of effective order $p$ (taking $\pi_{h}=\mathrm{id}$ ) or higher, but the converse is not true in general.

The analysis of the order conditions of the method $\hat{\psi}_{h}$ shows that many of them can be satisfied by $\pi_{h}$, so that $\psi_{h}$ must fulfill a much reduced set of restrictions [2,3]. In particular, if one takes a composition (5) for $\psi_{h}$, the number and complexity of the conditions to be verified by the coefficients $a_{i}, b_{i}$ is notably reduced. As a consequence, by considering both the kernel $\psi_{h}$ and the post-processor $\pi_{h}$ as compositions of basic integrators, highly efficient processed methods have been proposed [23,2,1,17]. Nevertheless, when $\pi_{h}$ is constructed as a composition (5), its computational cost is usually higher than that of $\psi_{h}$, and thus the use of the resulting processed scheme is restricted to situations where intermediate results are not frequently required. Otherwise the overall efficiency of the method is deteriorated.

To overcome this difficulty, in [3] a technique has been developed for obtaining approximations to the post-processor at virtually cost free and without loss of accuracy. The clue is to replace $\pi_{h}$ by a new integrator $\tilde{\pi}_{h} \simeq \pi_{h}$ obtained from the intermediate stages in the computation of $\psi_{h}$. As a result of the analysis carried out in [3], it is generally recommended to have a very accurate pre-processor $\pi_{h}^{-1}$ but, on the contrary, $\pi_{h}$ can safely be replaced by $\tilde{\pi}_{h}$, since the error introduced by the cheap approximation $\tilde{\pi}_{h}$ is of a purely local character and is not propagated along the evolution (contrarily to the error in $\pi_{h}^{-1}$ ).

Here we also address the following question: do Theorems 1 and 2 also hold for a composition $\psi_{h}$ of effective order $p \geqslant 3$ ? Observe that, in principle, Theorem 2 applies to the whole composition $\hat{\psi}_{h}$ of (6), but it would nonetheless be advantageous to have the negative coefficients restricted only to the composition $\pi_{h}$. For in that case the integration starts by computing $\pi_{h}^{-1}$ (which only involves positive coefficients), then $\psi_{h}$ (involving only positive coefficients) is evaluated once per step and finally an appropriate approximation to $\pi_{h}$ may be considered when output is required (even the crudest approximation $\pi_{h}=\mathrm{id}$ [13]). In this way the algorithm only involves stepping forward in time and could be applied even to PDEs with unbounded operators. The answer to the question posed before could also be useful in the search of efficient methods of order higher than 2 for systems that evolve in a semigroup, such as the heat equation in two space dimensions [17].

Also in quantum statistical mechanics, the partition function requires computing $Z=\operatorname{Tr}\left(\mathrm{e}^{-\beta H}\right)$, where $H$ is the Hamiltonian operator and $\beta$ is the inverse temperature [22]. In numerical calculations a processed composition algorithm may be used to approximate $\mathrm{e}^{-\beta H}$, and since the trace is invariant under similarity transformations, only the kernel is necessary to determine $Z$. If it involves only positive coefficients then it would be possible to build up higher order convergent algorithms for this class of problems.

In Section 3 we prove explicitly that this is not the case, so that any composition method (5) of effective order $p \geqslant 3$ contains necessarily some negative coefficients.

## 2. An elementary proof of Theorem 2

It is well known that, for each infinitely differentiable map $g: \mathbb{R}^{D} \rightarrow \mathbb{R}, g\left(\varphi_{h}(x)\right)$ admits an expansion of the form

$$
g\left(\varphi_{h}(x)\right)=g(x)+\sum_{k \geqslant 1} \frac{h^{k}}{k!} F^{k}[g](x)=\mathrm{e}^{h F}[g](x), \quad x \in \mathbb{R}^{D}
$$

where $F$ is the vector field (2). Similarly, for the map $\psi_{h}$ given in (5) one has

$$
g\left(\psi_{h}(x)\right)=\Psi_{h}[g](x)
$$

where [10]

$$
\begin{equation*}
\Psi_{h}=\exp \left(h a_{1} F_{A}\right) \exp \left(h b_{1} F_{B}\right) \cdots \exp \left(h a_{m} F_{A}\right) \exp \left(h b_{m} F_{B}\right) \tag{8}
\end{equation*}
$$

By repeated application of the Baker-Campbell-Hausdorff ( BCH ) formula to (8) we can obtain a series of differential operators $F_{h}=\sum_{k \geqslant 1} h^{k} F_{k}$ such that $\Psi_{h}=\exp \left(F_{h}\right)$, i.e., $\psi_{h}$ is formally the exact 1-flow of the vector field $F_{h}$. The scheme $\psi_{h}$ is of order $p \geqslant 3$ if $F_{1}=F_{A}+F_{B}$ and $F_{2}=F_{3}=0$. In terms of the coefficients $a_{i}, b_{i}$, this corresponds to the following order conditions:
order 1: $\quad \sum_{i=1}^{m} a_{i}=1, \quad \sum_{i=1}^{m} b_{i}=1 ;$
order $2: \quad \sum_{i=1}^{m} b_{i}\left(\sum_{j=1}^{i} a_{j}\right)=\frac{1}{2} ;$
order 3: $\quad \sum_{i=1}^{m-1} b_{i}\left(\sum_{j=i+1}^{m} a_{j}\right)^{2}=\frac{1}{3}, \quad \sum_{i=1}^{m} a_{i}\left(\sum_{j=i}^{m} b_{j}\right)^{2}=\frac{1}{3}$.
The proof of Theorems 1 and 2 provided by [20,22,9] is based precisely on the fact that a scheme of the form (5) with $m$ any finite positive integer and all the coefficients $a_{i}, b_{i}$ being positive cannot satisfy Eqs. (9).

In the particular case of the first order method $\chi_{h}=\varphi_{h}^{[B]} \circ \varphi_{h}^{[A]}$, the corresponding operator (8) is given by

$$
\begin{equation*}
\exp \left(h F_{A}\right) \exp \left(h F_{B}\right)=\exp \left(X_{h}\right)=\exp \left(h X_{1}+h^{2} X_{2}+h^{3} X_{3}+\cdots\right) \tag{10}
\end{equation*}
$$

with $X_{1}=F_{A}+F_{B}, X_{2}=\frac{1}{2}\left[F_{A}, F_{B}\right]$, etc., whereas for its adjoint scheme $\chi_{h}^{*}=\chi_{-h}^{-1}=\varphi_{h}^{[A]} \circ \varphi_{h}^{[B]}$ one has $g\left(\chi_{h}^{*}(x)\right)=\mathrm{e}^{-X_{-h}}[g](x)$. Here $\left[F_{A}, F_{B}\right]$ stands for the Lie bracket $F_{A} F_{B}-F_{B} F_{A}$. Since our aim is to get results valid for all pairs $F_{A}, F_{B}$ of arbitrary vector fields, then we must assume that the only linear dependencies among nested Lie brackets of $F_{A}$ and $F_{B}$ can be derived from the skew-symmetry and the Jacobi identity of the Lie brackets. In other words, $X_{k}, k \geqslant 1$, is an element of the graded free Lie algebra generated by the symbolic vector fields $F_{A}, F_{B}$, where both have degree one [18].

The crucial observation that leads to an alternative, elementary proof of Theorem 2 is the close connection existing between the splitting method (5)

$$
\begin{equation*}
\psi_{h}=\varphi_{b_{m} h}^{[B]} \circ \varphi_{a_{m} h}^{[A]} \circ \varphi_{b_{m-1} h}^{[B]} \circ \cdots \circ \varphi_{a_{2} h}^{[A]} \circ \varphi_{b_{1} h}^{[B]} \circ \varphi_{a_{1} h}^{[A]} \tag{11}
\end{equation*}
$$

and the composition of the first order method $\chi_{h}=\varphi_{h}^{[B]} \circ \varphi_{h}^{[A]}$ and its adjoint $\chi_{h}^{*}=\varphi_{h}^{[A]} \circ \varphi_{h}^{[B]}$ with different time steps [15]:

$$
\begin{equation*}
\psi_{h}=\chi_{\beta_{2 m} h}^{*} \circ \chi_{\beta_{2 m-1} h} \circ \cdots \circ \chi_{\beta_{2} h}^{*} \circ \chi_{\beta_{1} h} \circ \chi_{\beta_{0} h}^{*} \tag{12}
\end{equation*}
$$

Indeed, by inserting the explicit form of $\chi_{\beta_{i} h}$ and $\chi_{\beta_{i} h}^{*}$ in (12) we have

$$
\begin{aligned}
\psi_{h} & =\left(\varphi_{\beta_{2 m} h}^{[A]} \circ \varphi_{\beta_{2 m} h}^{[B]}\right) \circ\left(\varphi_{\beta_{2 m-1} h}^{[B]} \circ \varphi_{\beta_{2 m-1} h}^{[A]}\right) \circ \cdots \circ\left(\varphi_{\beta_{2} h}^{[A]} \circ \varphi_{\beta_{2} h}^{[B]}\right) \circ\left(\varphi_{\beta_{1} h}^{[B]} \circ \varphi_{\beta_{1} h}^{[A]}\right) \circ\left(\varphi_{\beta_{0} h}^{[A]} \circ \varphi_{\beta_{0} h}^{[B]}\right) \\
& =\varphi_{\beta_{2 m} h}^{[A]} \circ \varphi_{\left(\beta_{2 m}+\beta_{2 m-1}\right) h}^{[B]} \circ \varphi_{\left(\beta_{2 m-1}+\beta_{2 m-2}\right) h}^{[A]} \circ \cdots \circ \varphi_{\left(\beta_{2}+\beta_{1}\right) h}^{[B]} \circ \varphi_{\left(\beta_{1}+\beta_{0}\right) h}^{[A]} \circ \varphi_{\beta_{0} h}^{[B]},
\end{aligned}
$$

where in the last equality we have used the group property of the exact flow. If we put $\beta_{0}=\beta_{2 m}=0$ we recover expression (11) as soon as

$$
\begin{equation*}
a_{i}=\beta_{2 i-1}+\beta_{2 i-2}, \quad b_{i}=\beta_{2 i}+\beta_{2 i-1}, \quad i=1, \ldots, m \tag{13}
\end{equation*}
$$

Then

$$
\begin{equation*}
\sum_{i=1}^{m} a_{i}=\sum_{i=0}^{2 m} \beta_{i}=\sum_{i=1}^{m} b_{i} \tag{14}
\end{equation*}
$$

In consequence, composition (11) can be rewritten as (12) only if (14) holds. Consistency of both schemes require in fact that $\sum_{i=1}^{m} a_{i}=\sum_{i=1}^{m} b_{i}=\sum_{i=0}^{2 m} \beta_{i}=1$ and it has been shown in [15] that the order conditions for the coefficients $a_{i}, b_{i}$ to get a method of order $p \geqslant 1$ are equivalent to the order conditions for the $\beta_{i}$. In that case the operator $\Psi_{h}$ given in (8) can also be expressed as

$$
\begin{equation*}
\Psi_{h}=\exp \left(-X_{-\beta_{0} h}\right) \exp \left(X_{\beta_{1} h}\right) \exp \left(-X_{-\beta_{2} h}\right) \cdots \exp \left(X_{\beta_{2 m-1} h}\right) \exp \left(-X_{-\beta_{2 m} h}\right) \tag{15}
\end{equation*}
$$

and repeated application of the BCH formula gives

$$
\Psi_{h}=\exp \left(h f_{1,1} X_{1}+h^{2} f_{2,1} X_{2}+h^{3}\left(f_{3,1} X_{3}+f_{3,2}\left[X_{1}, X_{2}\right]\right)+\mathcal{O}\left(h^{4}\right)\right)
$$

where the coefficients $f_{k, j}$ are homogeneous polynomials of degree $k$ in the variables $\beta_{i}$. In particular we have

$$
\begin{equation*}
f_{1,1}=\sum_{i=0}^{2 m} \beta_{i}, \quad f_{2,1}=\sum_{i=0}^{2 m}(-1)^{i+1} \beta_{i}^{2}, \quad f_{3,1}=\sum_{i=0}^{2 m} \beta_{i}^{3} \tag{16}
\end{equation*}
$$

Conditions $f_{1,1}=1$ and $f_{n, j}=0$ for all $n \leqslant p$ are then sufficient for the method to be of order $p$.
Proof of Theorem 2. From the preceding discussion, it is clear that

$$
\begin{equation*}
f_{3,1}=\sum_{i=0}^{2 m} \beta_{i}^{3}=0 \tag{17}
\end{equation*}
$$

is a necessary condition to be satisfied by any method of order $p \geqslant 3$. We suppose that more than two $\beta_{i}$ are different from zero because $\beta_{1}^{3}+\beta_{2}^{3}=0$ together with the consistency condition $\beta_{1}+\beta_{2}=1$ have no real solution. Now (17) can be written as

$$
\sum_{i=1}^{m}\left(\beta_{2 i-1}^{3}+\beta_{2 i-2}^{3}\right)+\beta_{2 m}^{3}=\sum_{i=1}^{m}\left(\beta_{2 i-1}^{3}+\beta_{2 i-2}^{3}\right)=0
$$

for any positive integer $m$. In consequence $\beta_{2 j-1}^{3}+\beta_{2 j-2}^{3}$ has to be negative for some $1 \leqslant j \leqslant m$. But it is easy to verify that $\operatorname{sign}\left(x^{3}+y^{3}\right)=\operatorname{sign}(x+y)$ for any $x, y \in \mathbb{R}$, so that

$$
\begin{equation*}
a_{j}=\beta_{2 j-1}+\beta_{2 j-2}<0 \tag{18}
\end{equation*}
$$

for some $j$ such that $1 \leqslant j \leqslant m$. Similarly, we can write (17) as

$$
\beta_{0}^{3}+\sum_{i=1}^{m}\left(\beta_{2 i}^{3}+\beta_{2 i-1}^{3}\right)=\sum_{i=1}^{m}\left(\beta_{2 i}^{3}+\beta_{2 i-1}^{3}\right)=0
$$

so that $\beta_{2 k}^{3}+\beta_{2 k-1}^{3}<0$ for some $1 \leqslant k \leqslant m$, and again

$$
\begin{equation*}
b_{k}=\beta_{2 k}+\beta_{2 k-1}<0 \tag{19}
\end{equation*}
$$

## Distribution of the coefficients

We can get more information about the distribution of the negative coefficients in the composition (5) by applying a slightly more involved argument which, in fact, also provides another demonstration of Theorem 2.

If we denote $\alpha_{2 i-1}=a_{i}, \alpha_{2 i}=b_{i}, i=1, \ldots, m$, in the composition (5)

$$
\psi_{h}=\varphi_{\alpha_{2 m} h}^{[B]} \circ \varphi_{\alpha_{2 m-1} h}^{[A]} \circ \cdots \circ \varphi_{\alpha_{2} h}^{[B]} \circ \varphi_{\alpha_{1} h}^{[A]}
$$

then (13) implies

$$
\begin{equation*}
\alpha_{i}=\beta_{i}+\beta_{i-1}, \quad i=1, \ldots, 2 m \tag{20}
\end{equation*}
$$

where $\beta_{j}$ are the coefficients appearing in (12) $\left(\beta_{0}=\beta_{2 m}=0\right)$. Now all we need to prove Theorem 2 is to analyse how Eqs. (17) and (20) imply that at least one odd as well as one even $\alpha_{i}$ coefficients are negative.

As before, we assume that there are more than two nonvanishing coefficients $\beta_{i}$ and at least one of them is negative.
(a) Let us suppose first that only one coefficient is actually negative, say $\beta_{j}$, for some $0<j<2 m$. Then, from Eq. (17),

$$
\beta_{j}=-\left(\sum_{i \neq j} \beta_{i}^{3}\right)^{1 / 3}
$$

so that $\left|\beta_{j}\right|>\beta_{i}$ for all $i \neq j$. Therefore

$$
\alpha_{j}=\beta_{j}+\beta_{j-1}<0 \quad \text { and } \quad \alpha_{j+1}=\beta_{j+1}+\beta_{j}<0
$$

i.e., two consecutive $\alpha_{k}$ coefficients are negative, and thus at least one $a_{j}$ and one $b_{j}$ are negative.
(b) Suppose now that the negative coefficients are $\beta_{j_{1}}, \beta_{j_{2}}, \ldots, \beta_{j_{k}}$ with $j_{1}<j_{2}<\cdots<j_{k}$.
(b.1) If

$$
\begin{equation*}
\beta_{j_{i}-1}<\left|\beta_{j_{i}}\right| \quad \text { and } \quad\left|\beta_{j_{i}}\right|>\beta_{j_{i}+1} \tag{21}
\end{equation*}
$$

for some $j_{i} \in \mathcal{I} \equiv\left\{j_{1}, j_{2}, \ldots, j_{k}\right\}$ then also two consecutive coefficients $\alpha_{k}$ are negative, namely $\alpha_{j_{i}}$ and $\alpha_{j_{i}+1}$.
(b.2) On the other hand, when (21) does not hold for any $j_{i} \in \mathcal{I}$, then the following situations are possible:
(i) if $\beta_{j_{i}-1}<\left|\beta_{j_{i}}\right|$, then $\left|\beta_{j_{i}}\right|<\beta_{j_{i}+1}$;
(ii) if $\left|\beta_{j_{i}}\right|>\beta_{j_{i}+1}$, then $\beta_{j_{i}-1}>\left|\beta_{j_{i}}\right|$;
(iii) finally, $\beta_{j_{i}-1}>\left|\beta_{j_{i}}\right|$ and $\beta_{j_{i}+1}>\left|\beta_{j_{i}}\right|$.

Let us suppose that $\beta_{j_{i}+1} \neq \beta_{j_{i+1}-1}$ for all $j_{i}$. Then

$$
\begin{aligned}
\beta_{j_{k}}= & -\left(\left(\beta_{j_{1}-1}^{3}+\beta_{j_{1}}^{3}+\beta_{j_{1}+1}^{3}\right)+\left(\beta_{j_{2}-1}^{3}+\beta_{j_{2}}^{3}+\beta_{j_{2}+1}^{3}\right)+\cdots\right. \\
& \left.+\left(\beta_{j_{k-1}-1}^{3}+\beta_{j_{k-1}}^{3}+\beta_{j_{k-1}+1}^{3}\right)+\sum^{\prime} \beta_{i}^{3}\right)^{1 / 3}
\end{aligned}
$$

where $\sum^{\prime}$ contains the remaining terms (including $\beta_{j_{k}-1}$ and $\beta_{j_{k}+1}$ ). Since $\beta_{j_{l}-1}^{3}+\beta_{j_{l}}^{3}+\beta_{j_{l}+1}^{3}>$ 0 for $l=1, \ldots, k-1$, then clearly

$$
\left|\beta_{j_{k}}\right|>\beta_{j_{k}-1} \quad \text { and } \quad\left|\beta_{j_{k}}\right|>\beta_{j_{k}+1}
$$

in contradiction with hypothesis (b.2). Therefore $\beta_{j_{i}+1}=\beta_{j_{i+1}-1}$ for some $j_{i}$. Let us suppose, without loss of generality, that they correspond to the first $l+1$ coefficients $\beta_{j_{i}}$. Then, the only
possible sequence (different from those considered before) has to be $\beta_{j_{i}-1}, \beta_{j_{i}}, \beta_{j_{i}+1}, \beta_{j_{i}+2}$, $\beta_{j_{i}+3}, \ldots, \beta_{j_{i}+2 l-1}, \beta_{j_{i}+2 l}, \beta_{j_{i}+2 l+1}$ such that

$$
\begin{aligned}
& \beta_{j_{i}-1}<\left|\beta_{j_{i}}\right|<\beta_{j_{i}+1}, \\
& \left|\beta_{j_{i}+2}\right|<\beta_{j_{i}+1}, \quad\left|\beta_{j_{i}+2}\right|<\beta_{j_{i}+3}, \quad \ldots \\
& \beta_{j_{i}+2 l-1}>\left|\beta_{j_{i}+2 l}\right|>\beta_{j_{i}+2 l+1}
\end{aligned}
$$

where $\beta_{j_{i}+2 k}, k=0,1, \ldots, l$ are the negative coefficients. Then

$$
\alpha_{j_{i}}=\beta_{j_{i}}+\beta_{j_{i}-1}<0, \quad \alpha_{j_{i}+2 l+1}=\beta_{j_{i}+2 l+1}+\beta_{j_{i}+2 l}<0
$$

Also in this case at least one $a_{i}$ and one $b_{i}$ are negative because $j_{i}$ and $j_{i}+2 l+1$ differ in an odd number.

Notice that this is the only situation where two negative $\alpha_{i}$ coefficients in a given method do not stay in consecutive places. We have checked several composition methods published in the literature having observed that this occurrence is in fact quite rare: it is very much frequent that at least two consecutive $\alpha_{i}$ coefficients are negative, and this discussion provides an explanation of the phenomenon.

## 3. Compositions of effective order $p \geqslant 3$

As with the composition $\psi_{h}$ in (5), let us consider a post-processor $\pi_{h}$ in (6) formally as the exact 1-flow of the vector field $P_{h}$, i.e., $g\left(\pi_{h}(x)\right)=\mathrm{e}^{P_{h}}[g](x)$ for all $g$, with $P_{h}=\sum_{k \geqslant 1} h^{k} P_{k}$. Then one has for the processed scheme $g \circ \hat{\psi}_{h}=\widehat{\Psi}_{h}[g]$, where $\widehat{\Psi}_{h}=\exp \left(\widehat{F}_{h}\right)$ and the vector field $\widehat{F}_{h}$ can be determined from the relation

$$
\begin{equation*}
\exp \left(\widehat{F}_{h}\right)=\exp \left(-P_{h}\right) \exp \left(F_{h}\right) \exp \left(P_{h}\right) \tag{22}
\end{equation*}
$$

With respect to the vector field $P_{h}$, it is natural to choose it as an element of the graded free Lie algebra generated by $F_{A}$ and $F_{B}$. Thus, up to order two in $h$,

$$
\begin{equation*}
P_{h}=h\left(c_{1} F_{A}+c_{2} F_{B}\right)+h^{2} c_{3} X_{2}+\mathcal{O}\left(h^{3}\right) \tag{23}
\end{equation*}
$$

with $c_{i}$ free parameters. Notice that $\pi_{h}$ can be approximated by a composition (5), or equivalently, $\exp \left(P_{h}\right)$ by the product (8). However, if $c_{1} \neq c_{2}$ then $\sum_{i} a_{i}=c_{1} \neq c_{2}=\sum_{i} b_{i}$ and composition (12) cannot be used (as is the case, for instance, of the Strang splitting (7)). On the other hand, since $c_{1} F_{A}+c_{2} F_{B}=\left(c_{2}-c_{1}\right) F_{B}+c_{1} X_{1}=\left(c_{1}-c_{2}\right) F_{A}+c_{2} X_{1}$, from (23) it is possible to write

$$
\begin{equation*}
\mathrm{e}^{P_{h}}=\mathrm{e}^{h c_{1} X_{1}} \mathrm{e}^{h c F_{B}} \mathrm{e}^{h^{2} d_{1} X_{2}}+\mathcal{O}\left(h^{3}\right) \tag{24}
\end{equation*}
$$

as well as

$$
\begin{equation*}
\mathrm{e}^{P_{h}}=\mathrm{e}^{h c_{2} X_{1}} \mathrm{e}^{-h c F_{A}} \mathrm{e}^{h^{2} d_{2} X_{2}}+\mathcal{O}\left(h^{3}\right) \tag{25}
\end{equation*}
$$

where $c=c_{2}-c_{1}$ and $d_{1}, d_{2}$ are parameters depending on $c_{1}, c_{2}, c_{3}$. Since the processed scheme $\hat{\psi}_{h}$ is of conventional order $p$, then $\widehat{\Psi}_{h}=\exp \left(h X_{1}\right)+\mathcal{O}\left(h^{p+1}\right)$ and $\mathrm{e}^{h X_{1}} \widehat{\Psi}_{h}=\widehat{\Psi}_{h} \mathrm{e}^{h X_{1}}+\mathcal{O}\left(h^{p+1}\right)$, so that $\exp \left(h c_{i} X_{1}\right)$ in (24) and (25) can be safely removed without loss of generality and thus we take $c_{1}=0$ or $c_{2}=0$.

Now we are in disposition to establish and prove the main result of this section.

Theorem 3. At least one of the $a_{i}$ as well as one of the $b_{i}$ coefficients have to be negative in the composition (5) if $\psi_{h}$ is the kernel of a processed method of order (or equivalently if $\psi_{h}$ is of effective order) $p \geqslant 3$.

Proof. Let $\psi_{h}$ be a composition (5) of effective order 3 with, say, all $a_{i}$ positive. Then it is possible to construct a vector field $P_{h}$ such that (24) (with $c_{1}=0$ ) holds and therefore (22) leads to

$$
\begin{equation*}
\mathrm{e}^{-h^{2} d_{1} X_{2}} \mathrm{e}^{-h c F_{B}} \Psi_{h} \mathrm{e}^{h c F_{B}} \mathrm{e}^{h^{2} d_{1} X_{2}}=\mathrm{e}^{\widehat{F_{h}}}=\mathrm{e}^{h\left(F_{A}+F_{B}\right)}+\mathcal{O}\left(h^{4}\right) \tag{26}
\end{equation*}
$$

or equivalently

$$
\begin{align*}
\bar{\Psi}_{h} & \equiv \mathrm{e}^{-h c F_{B}} \Psi_{h} \mathrm{e}^{h c F_{B}}=\mathrm{e}^{h^{2} d_{1} X_{2}} \mathrm{e}^{h\left(F_{A}+F_{B}\right)} \mathrm{e}^{-h^{2} d_{1} X_{2}}+\mathcal{O}\left(h^{4}\right) \\
& =\exp \left(h\left(F_{A}+F_{B}\right)-h^{3} d_{1}\left[X_{1}, X_{2}\right]\right)+\mathcal{O}\left(h^{4}\right), \tag{27}
\end{align*}
$$

where $\Psi_{h}$ is given by (8). Notice that all coefficients $a_{i}$ in $\bar{\Psi}_{h}$ are positive, since $\bar{\Psi}_{h}$ is associated with the composition map

$$
\begin{equation*}
\bar{\psi}_{h}=\varphi_{c h}^{[B]} \circ \varphi_{b_{m} h}^{[B]} \circ \varphi_{a_{m} h}^{[A]} \circ \cdots \circ \varphi_{b_{1} h}^{[B]} \circ \varphi_{a_{1} h}^{[A]} \circ \varphi_{-c h}^{[B]}, \tag{28}
\end{equation*}
$$

which, as previously, can be written as a composition of the first order method $\chi_{h}$ and its adjoint $\chi_{h}^{*}$ with coefficients $\bar{\beta}_{i}$ :

$$
\begin{equation*}
\bar{\psi}_{h}=\chi_{\bar{\beta}_{2 k} h}^{*} \circ \chi_{\bar{\beta}_{2 k-1} h} \circ \cdots \circ \chi_{\bar{\beta}_{2} h}^{*} \circ \chi_{\bar{\beta}_{1} h} \circ \chi_{\bar{\beta}_{0} h}^{*} \tag{29}
\end{equation*}
$$

with $\bar{\beta}_{0}=\bar{\beta}_{2 k}=0$. Since the coefficient of $X_{3}$ is zero in (27), it is clear that

$$
\begin{equation*}
\bar{f}_{3,1}:=\sum_{i=0}^{2 k} \bar{\beta}_{i}^{3}=0 \tag{30}
\end{equation*}
$$

But, as we know from the proof of Theorem 2 provided in Section 2, this condition cannot be satisfied with all coefficients $a_{i}$ positive.

Similarly, if we assume that all $b_{i}$ are positive then the same argument applied to the post-processor (25) leads to the same contradiction.

In fact, the explicit expression of the effective order conditions up to order 3 can be derived in the following way. By inserting (24) in (22) and applying the BCH formula one finds

$$
\begin{align*}
\widehat{\Psi}_{h}= & \exp \left(h X_{1}+h^{2}\left(f_{2,1}+2 c\right) X_{2}+h^{3}\left(\left(f_{3,1}+3 c\left(f_{2,1}+c\right)\right) X_{3}\right.\right. \\
& \left.\left.+\left(f_{3,2}-\frac{1}{2} c\left(f_{2,1}+c\right)+d_{1}\right)\left[X_{1}, X_{2}\right]\right)\right)+\mathcal{O}\left(h^{4}\right) \tag{31}
\end{align*}
$$

and a second order method is obtained by taking $c=-\frac{1}{2} f_{2,1}$. If we substitute this value in (31) and take $d_{1}$ such that the coefficient of [ $X_{1}, X_{2}$ ] vanishes, then

$$
\begin{equation*}
\widehat{\Psi}_{h}=\exp \left(h X_{1}+h^{3}\left(f_{3,1}-\frac{3}{4} f_{2,1}^{2}\right) X_{3}\right)+\mathcal{O}\left(h^{4}\right) \tag{32}
\end{equation*}
$$

This same result is obtained if one considers a post-processor such that (25) holds (with $c_{2}=0$ ) instead of (24). In summary, it is clear that

$$
\begin{equation*}
f_{3,1}-\frac{3}{4} f_{2,1}^{2}=0 \tag{33}
\end{equation*}
$$

is the only condition to be satisfied by a composition to be of effective order three. This condition is equivalent to the kernel condition at order three presented in [2] using a different basis of the Lie algebra.

Example. Let us consider the composition map

$$
\begin{equation*}
\psi_{h}=\varphi_{(1-b) h}^{[B]} \circ \varphi_{(1-a) h}^{[A]} \circ \varphi_{b h}^{[B]} \circ \varphi_{a h}^{[A]}=\chi_{\beta_{4} h}^{*} \circ \chi_{\beta_{3} h} \circ \chi_{\beta_{2} h}^{*} \circ \chi_{\beta_{1} h} \circ \chi_{\beta_{0} h}^{*}, \tag{34}
\end{equation*}
$$

with $\beta_{0}=0, \beta_{1}=a, \beta_{2}=b-a, \beta_{3}=1-b, \beta_{4}=0$, and the consistency conditions already imposed. This composition cannot be of order three (there are not enough parameters to solve all the order conditions (9)), yet it could be of effective order three if condition (33) is satisfied, which in this case reads

$$
1-12 a b(1-a)(1-b)=0
$$

But it turns out that this equation has no real solution if $a \in(0,1)$ as well as if $b \in(0,1)$.

## 4. Other classes of composition methods

The previous results can be generalized in different contexts. For instance, let us consider a partitioned scheme built up using finite linear combinations of splitting methods of the form (5), i.e.,

$$
\begin{equation*}
\psi_{h}=\sum_{k=1}^{K} \gamma_{k} \psi_{h, k} \tag{35}
\end{equation*}
$$

where

$$
\begin{equation*}
\psi_{h, k}=\varphi_{b_{k, h} h}^{[B]} \circ \varphi_{a_{k, h} h}^{[A]} \circ \cdots \circ \varphi_{b_{k, 1} h}^{[B]} \circ \varphi_{a_{k, 1} h}^{[A]} \tag{36}
\end{equation*}
$$

and it is assumed that $\sum_{k} \gamma_{k}=1$ and $\sum_{i} a_{k, i}=\sum_{i} b_{k, i}, k=1, \ldots, K$. The generalization provided by the following theorem establishes that even with partitioned schemes of the form (35) each basic flow in a convex partition ( $\gamma_{k}>0$ for every $1 \leqslant k \leqslant K$ ) must be applied for at least one backward fractional time step. On the other hand, simple polynomial extrapolation of the leapfrog method (7) shows that if $\gamma_{k}<0$ all the coefficients $a_{k, i}, b_{k, i}$ may indeed be positive.

Theorem 4 [9]. If $p$ and $K$ are positive integers such that $p \geqslant 3$ and $K \geqslant 2$, and $\gamma_{k}>0$ for $k=1, \ldots, K$, then at least one of the coefficients $a_{k, i}$ as well as one of the $b_{k, i}$ have to be negative in the composition (36) if $\psi_{h}$ has order, or effective order, $p \geqslant 3$.

Proof. Also in this case the proof is quite elementary. Since $\sum_{i} a_{k, i}=\sum_{i} b_{k, i}$ we can write

$$
\begin{equation*}
\psi_{h, k}=\chi_{\beta_{k, 2 h} h}^{*} \circ \chi_{\beta_{k, 2 n-1} h} \circ \cdots \circ \chi_{\beta_{k, 2} h}^{*} \circ \chi_{\beta_{k, 1} h} \circ \chi_{\beta_{k, 0} h}^{*} \tag{37}
\end{equation*}
$$

and by following the same procedure as previously we find that instead of (16) the necessary condition for (35) to be a method of order three or higher is now

$$
\sum_{k=1}^{K} \gamma_{k} \sum_{i=1}^{2 n} \beta_{k, i}^{3}=0 .
$$

Since $\gamma_{k}>0, k=1, \ldots, K$, there exists some $j, 1 \leqslant j \leqslant K$, such that

$$
\sum_{i=1}^{2 n} \beta_{j, i}^{3} \leqslant 0
$$

and the previous proof can be applied.
The two-term splitting analyzed so far can be seen as a special instance of a $k$-term splitting of the vector field $F, F=F_{1}+F_{2}+\cdots+F_{k}$. Suppose we have a scheme of the form

$$
\begin{align*}
\psi_{h}= & \varphi_{a_{m, k} h}^{[k]} \circ \cdots \circ \varphi_{a_{m, 2} h}^{[2]} \circ \varphi_{a_{m, 1} h}^{[1]} \circ \cdots \circ \varphi_{a_{2, k} h}^{[k]} \circ \cdots \circ \varphi_{a_{2,2} h}^{[2]} \circ \varphi_{a_{2,1} h}^{[1]} \\
& \circ \varphi_{a_{1, k} h}^{[k]} \circ \cdots \circ \varphi_{a_{1,2} h}^{[2]} \circ \varphi_{a_{1,1} h}^{[1]}, \tag{38}
\end{align*}
$$

where $\varphi_{h}^{[l]}$ stands for the exact $h$-flow of the vector field $F_{l}$. If the composition (38) is of order, or effective order, $p \geqslant 3$ for all choices of operators $F_{1}, \ldots, F_{k}$, then clearly

$$
\min _{1 \leqslant i \leqslant m} a_{i, l}<0, \quad l=1, \ldots, k
$$

Consider now composition maps of the Strang splitting $\psi_{h, 2}$ given by (4) (or with the roles of the flows $\varphi_{h}^{[A]}$ and $\varphi_{h}^{[B]}$ interchanged), i.e.,

$$
\begin{equation*}
\psi_{h}=\psi_{\beta_{m} h, 2} \circ \psi_{\beta_{m-1} h, 2} \circ \cdots \circ \psi_{\beta_{1} h, 2} \circ \psi_{\beta_{0} h, 2} \tag{39}
\end{equation*}
$$

The series of differential operators $S_{h}$ associated with the integrator $\psi_{h, 2}$, i.e., such that $g \circ \psi_{h, 2}=S_{h}[g]$, can be written as $S_{h}=\exp \left(X_{h}\right)$, where $X_{h}=h X_{1}+h^{3} X_{3}+h^{5} X_{5}+\cdots, X_{1}=F$ and

$$
\Psi_{h}=\exp \left(X_{\beta_{0} h}\right) \exp \left(X_{\beta_{1} h}\right) \cdots \exp \left(X_{\beta_{m-1} h}\right) \exp \left(X_{\beta_{m} h}\right)
$$

Now, by repeated application of the BCH formula,

$$
\Psi_{h}=\exp \left(h f_{1,1} X_{1}+h^{3} f_{3,1} X_{3}+\mathcal{O}\left(h^{4}\right)\right)
$$

where

$$
f_{1,1}=\sum_{i=0}^{m} \beta_{i}, \quad f_{3,1}=\sum_{i=0}^{m} \beta_{i}^{3}
$$

Therefore $f_{1,1}=1, f_{3,1}=0$ are necessary conditions for $\psi_{h}$ to be of order $p \geqslant 3$. In fact, since $f_{2,1}=0$ in this case, they are also the conditions to be satisfied by $\psi_{h}$ to have effective order $p=3$ and the following theorem can be established.

Theorem 5. If $p$ is any positive integer such that $p \geqslant 3$ and $\psi_{h, 2}$ is the Strang splitting (or Störmer/Verlet scheme) (4), then at least two consecutive coefficients $a_{i}, b_{i}$ have to be negative in the composition (39) (when it is expressed in terms of the basic flows $\varphi_{h}^{[A]}, \varphi_{h}^{[B]}$ ) if $\psi_{h}$ is of order, or effective order, $p$. Even more, at least two coefficients $a_{i_{1}}, a_{i_{2}}$ have to be negative.

Proof. By substituting in (39) the expression of the basic method $\psi_{\beta_{i} h, 2}=\varphi_{\beta_{i} h / 2}^{[A]} \circ \varphi_{\beta_{i} h}^{[B]} \circ \varphi_{\beta_{i} h / 2}^{[A]}$, we obtain a composition of the type (5) with

$$
\begin{equation*}
b_{i}=\beta_{i}, \quad a_{i}=\frac{1}{2}\left(\beta_{i}+\beta_{i-1}\right), \quad i=1, \ldots, m \tag{40}
\end{equation*}
$$

if $\beta_{0}=\beta_{m}=0$. It is immediate to check that if there exists one negative coefficient, say $\beta_{j}<0,1 \leqslant j \leqslant$ $m-1$, and

$$
\begin{equation*}
\left|\beta_{j}\right|>\beta_{j-1} \tag{41}
\end{equation*}
$$

then $a_{j}<0, b_{j}<0$, whereas if

$$
\begin{equation*}
\left|\beta_{j}\right|>\beta_{j+1} \tag{42}
\end{equation*}
$$

then $b_{j}<0, a_{j+1}<0$. In other words, as soon as one of the $\beta_{i}$ is negative and its absolute value is higher than the previous one or the next one then the corresponding composition (5) has, at least, two consecutive coefficients which are negative.

Let us analyse the different possibilities arising from the order condition $f_{3,1}=0$.
(i) First, suppose there exists only one negative coefficient $\beta_{j}<0,1 \leqslant j \leqslant m-1$. Then

$$
\beta_{j}=-\left(\sum_{i \neq j} \beta_{i}^{3}\right)^{1 / 3}
$$

and both conditions (41) and (42) are satisfied so that, according to the previous discussion, $a_{j}<0$, $b_{j}<0$ and $a_{j+1}<0$.
(ii) Suppose now that there are $k \geqslant 2$ negative coefficients, $\beta_{j_{1}}, \beta_{j_{2}}, \ldots, \beta_{j_{k}}<0$ such that they do not satisfy conditions (41) and (42). Observe that they cannot be consecutive, otherwise either (41) or (42) are satisfied. Then we can write condition $f_{3,1}=0$ as

$$
\beta_{j_{k}}=-\left(\left(\beta_{j_{1}-1}^{3}+\beta_{j_{1}}^{3}\right)+\cdots+\left(\beta_{j_{k-1}-1}^{3}+\beta_{j_{k-1}}^{3}\right)+\sum^{\prime} \beta_{i}^{3}\right)^{1 / 3}
$$

where $\sum^{\prime}$ contains the remaining terms, including $\beta_{j_{k}-1}$ and $\beta_{j_{k}+1}$. Since $\beta_{j_{i}-1}^{3}+\beta_{j_{i}}^{3}>0, i=$ $1, \ldots, k-1$, then $\beta_{j_{k}-1}<\left|\beta_{j_{k}}\right|, \beta_{j_{k}+1}<\left|\beta_{j_{k}}\right|$, conditions (41) and (42) are in fact satisfied by $\beta_{j_{k}}$ and therefore $b_{j_{1}}, \ldots, b_{j_{k}-1}, a_{j_{k}}, b_{j_{k}}, a_{j_{k}+1}<0$.
(iii) Finally, consider the case in which $\beta_{j_{1}}, \beta_{j_{2}}, \ldots, \beta_{j_{k}}<0(k \geqslant 2)$ and only one of the coefficients $\beta_{j_{i}}, i=1, \ldots, k$, satisfies either (41) or (42). For instance, suppose that $\beta_{j_{1}}$ is such that $\left|\beta_{j_{1}}\right|>\beta_{j_{1}-1}$ (and therefore $a_{j_{1}}<0$ ). Then, condition $f_{3,1}=0$ can be expressed as

$$
\left(\beta_{j_{1}}^{3}+\beta_{j_{1}+1}^{3}\right)+\sum_{i \neq j_{1}, j_{1}+1} \beta_{i}^{3}=0
$$

but $\beta_{j_{1}}^{3}+\beta_{j_{1}+1}^{3}>0$, since (42) is not satisfied by $\beta_{j_{1}}$, so that

$$
\left(\beta_{j_{2}}^{3}+\beta_{j_{2}+1}^{3}\right)+\cdots+\left(\beta_{j_{k}}^{3}+\beta_{j_{k}+1}^{3}\right)+\sum^{\prime} \beta_{i}^{3}<0
$$

where, as before, $\sum^{\prime}$ contains the remaining (positive) terms. In consequence, there must exist some $2 \leqslant i \leqslant k$ such that $\beta_{j_{i}}^{3}+\beta_{j_{i}+1}^{3}<0$. Therefore we have at least $b_{j_{1}}, \ldots, b_{j_{k}}<0$ and $a_{j_{1}}, a_{j_{i}+1}<0$.

If $\beta_{j_{1}}$ satisfies (42) instead, a similar strategy applies and the same conclusion follows.

This result, together with the discussion of Section 2 justifies why it is so frequent that at least two consecutive coefficients are negative.

## 5. Composition methods with all coefficients being positive

Let us consider now the second order differential equation

$$
\begin{equation*}
y^{\prime \prime}=g(y) \tag{43}
\end{equation*}
$$

which can be written in the form (1) by taking $x \equiv\left(x_{1}, x_{2}\right)=\left(y, y^{\prime}\right)$ and $f_{A}(x)=\left(x_{2}, 0\right), f_{B}(x)=$ $\left(0, g\left(x_{1}\right)\right)$, or equivalently

$$
F_{A} \equiv x_{2} \frac{\partial}{\partial x_{1}}, \quad F_{B} \equiv g\left(x_{1}\right) \frac{\partial}{\partial x_{2}}
$$

This equation frequently appears in relevant problems arising in classical and quantum mechanics: there the operator $F_{A}$ is related to the kinetic energy (quadratic in momenta) and $F_{B}$ is associated with the potential energy. Now the flow corresponding to $F_{C} \equiv\left[F_{B},\left[F_{A}, F_{B}\right]\right]$ is explicitly and exactly computable and, in addition, $\left[F_{B}, F_{C}\right]=0$, so that it makes sense to compute the 1-flow $\varphi_{b h, c h^{3}}^{[B, C]}$ associated with the vector field $h b F_{B}+c h^{3} F_{C}$ and include it into the composition (5):

$$
\begin{equation*}
\psi_{h}=\varphi_{b_{m} h, c_{m} h^{3}}^{[B, C]} \circ \varphi_{a_{m} h}^{[A]} \circ \cdots \circ \varphi_{b_{1} h, c_{1} h^{3}}^{[B, C]} \circ \varphi_{a_{1} h}^{[A]} . \tag{44}
\end{equation*}
$$

In this case $\psi_{h}$ cannot always be written as the composition of a first order scheme and its adjoint, and Theorem 2 does not necessarily apply. For instance

$$
\begin{equation*}
\psi_{h}=\varphi_{h / 6}^{[B]} \circ \varphi_{h / 2}^{[A]} \circ \varphi_{2 h / 3, h^{3} / 72}^{[B, C]} \circ \varphi_{h / 2}^{[A]} \circ \varphi_{h / 6}^{[B]} \tag{45}
\end{equation*}
$$

is a method of order four [11] and

$$
\begin{equation*}
\psi_{h}=\varphi_{h / 2}^{[A]} \circ \varphi_{h, h^{3} / 24}^{[B, C]} \circ \varphi_{h / 2}^{[A]} \tag{46}
\end{equation*}
$$

is a method of effective order four [19]. In the last case we can write $\psi_{h}=\chi_{h / 2}^{*} \circ \chi_{h / 2}$ with $\chi_{h} \equiv \varphi_{h, h^{3} / 6}^{[B, C]} \circ$ $\varphi_{h}^{[A]}$. However, if we analyse the corresponding operator $\exp \left(X_{h}\right)=\exp \left(h X_{1}+h^{2} X_{2}+h^{3} X_{3}+\cdots\right)$ associated with $\chi_{h}$, we find that $X_{3}=\left[X_{1}, X_{2}\right] / 6$. Then $X_{3}$ is not an independent element and its contribution can be cancelled with a proper choice of the map $\pi_{h}$, thus giving a fourth-order method.

Numerical experiments suggest that this is the highest order one can get with the composition (44) with positive coefficients and a rigorous proof is at present under investigation. However, methods of effective order six as well as of order six are known to exist with all coefficients $b_{i}$ being positive.

On the other hand, if we consider a Hamiltonian system of the form

$$
H=T(p)+V(q)
$$

with $T$ quadratic in $p$ and $V(q)$ a polynomial function up to degree four in $q$ (or, in general, if $g(y)$ is a polynomial function up to degree three), then $F_{E} \equiv\left[F_{A},\left[F_{A},\left[F_{A},\left[F_{A}, F_{B}\right]\right]\right]\right]$ vanish or depends only on the momenta, i.e., $\left[F_{A}, F_{E}\right]=0$, and its flow can be computed exactly. In addition $F_{D} \equiv\left[F_{B},\left[F_{B},\left[F_{A},\left[F_{A}, F_{B}\right]\right]\right]\right]$ depends only on the coordinates and thus $\left[F_{B}, F_{D}\right]=0$. Thus one may consider composition maps involving the 1 -flows $\varphi_{a h, e h^{5}}^{[A, E]}, \varphi_{b h, c h^{3}, d h^{5}}^{[B, C, D]}$ corresponding to the vector fields $h a F_{A}+e h^{5} F_{E}$ and $h b F_{B}+c h^{3} F_{C}+d h^{5} F_{D}$, respectively. In particular, the generalised leapfrog splitting scheme

$$
\begin{equation*}
\psi_{h}=\varphi_{h / 2, e h^{5}}^{[A, E]} \circ \varphi_{h, c h^{3}, d h^{5}}^{[B, C, D]} \circ \varphi_{h / 2, e h^{5}}^{[A, E]} \tag{47}
\end{equation*}
$$

with $c=\frac{1}{24}, d=\frac{1}{1440}, e=\frac{1}{2880}$ is a method of effective order six, since these coefficients satisfy the kernel conditions collected in [2] up to this order.

We should recall that methods (45)-(47) are particular examples of composition schemes involving only positive coefficients. The possible existence of other families of composition methods of order $p \geqslant 3$ with positive coefficients is, at the time being, an open question of great interest, for instance, in the numerical integration of nonreversible systems.

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Since the completion of this work, S.A. Chin has published a different proof of Theorem 3; see [5].

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